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WITH CERTAIN TYPES OF NOISE AND SYSTEMATIC
ERRORS COMMON TO MANY HADAMARD TRANSFORM
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FINAL REPORT

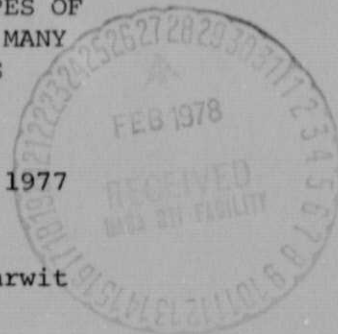
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PROCEDURES FOR DEALING WITH CERTAIN TYPES OF
NOISE AND SYSTEMATIC ERRORS COMMON TO MANY
HADAMARD TRANSFORM OPTICAL SYSTEMS

NASA Grant Langley NSG-1263

December 1, 1975 through November 30, 1977

Principal Investigator: Prof. Martin Harwit



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Summary

We have investigated procedures for dealing with a variety of different sources of noise and systematic errors that are commonly found in Hadamard transform and also in other multiplex systems. We were able to set up criteria for comparing instruments in the Fourier and Hadamard transform families through the use of a commonly applicable mathematical formalism. The results obtained in this study have led to the completion of one Ph.D. thesis by a Cornell University graduate student, Ming-Hing Tai, and we have also published four journal articles under this same contract. The most recent one of these which has been accepted for publication in 1978 appears here as Appendix. It contains most of the results we have been able to derive, in collaboration with Dr. N. J. A. Sloane, on sources of noise and errors, and the correction procedures that they call for.

1. Introduction

Our initial proposal submitted to NASA Langley in 1975 undertook to investigate a number of sources of noise and error correcting procedures characteristic of Hadamard transform optical systems. Specific questions to which we proposed to address ourselves involved the following four topics--as listed verbatim in the original proposal.

The questions of main concern are these:

1. A spectrum or picture can always be 'framed' within a region of constant known brightness. In spectroscopy, for example, the ends of a spectrum can be blocked off, so that m spectral elements on each end are known to have zero intensity. If N measurements are then made to determine the $N-2m$ unknown spectral intensities, there is a redundancy that enters because of the $2m$ constraints on the spectrum's end elements. Under these conditions, what is the optimum choice of the number m that will permit recovery of the best spectrum? Is there a better choice of an end pattern than a zero level constant signal?

2. Continuing in the same vein, we can ask to what extent the end elements can be used to help reduction of spectral noise due to noise spikes in the data. We mentioned above that a noise spike can be eliminated with good results and replaced by a mean data value. However, if the end elements are known to be zero intensity, it should be possible to eliminate spikes more judiciously by forcing the replaced readings to replace zeroes in the end positions of the spectrum.

3. Turning to errors in the mask construction, we need to know

i) What is the effect of random errors in the width of open and closed elements?

ii) What is the effect when the open elements, say, are systematically larger than the closed ones? This is an effect that can be brought about by excessive 'undercutting' in the photo-etch process frequently used in mask production.

iii) What errors are brought about when the steps in the mask motion systematically differ by a small amount from the size of the encoding element? In a typical application, the mask may have a dimension measuring some N elements. At the end of the pass, however, the mask has only been moved some $(N - \epsilon)$ elements, where, say, $\epsilon < 1/2$. What is the cumulative effect of this kind of error? How can it best be compensated? In practice stepper motors are driven by an accurate screw, and the errors discussed here are errors inherent in the construction of the screw.

4. In addition we also undertook "a study of a number of important limitations of this family of instruments that have not been attacked to date. Some of these correspond to related limitations found in Fourier transform instruments."

2. Results

We have been able to obtain an improved understanding of all of these effects, and have in fact been able to obtain analytical insights into a larger variety of sources of noise and error than had been originally anticipated.

We start with item 4 listed above. Here two publications have appeared in print (see section 3 below for publications 1 and 2). These two papers show the relative performance of Fourier and Hadamard transform spectrometers operated under identical detector-noise-limited conditions. Two main conclusions emerge from these papers. First, the singly-encoded Hadamard transform systems encode with an efficiency that is roughly $\sqrt{2}$ times higher than the corresponding efficiency obtained with a Michelson interferometric spectrometer. Second, the Michelson interferometer is less efficient than presented in standard texts. This second statement is not of the greatest importance because there had been some papers in the literature that had derived the correct performance figures; unfortunately, however, these papers had been widely ignored. Our contribution on this score was primarily to focus attention on the omission.

Among the topics listed under section 1, item 3., systematic means for dealing with mask defects of the type 3 i) and 3 iii) are treated exhaustively in publication 4 (see Appendix A and the list in section 3 below). In that paper written with Dr. Neil J. A. Sloane of the Bell Laboratories, Murray Hill, N.J., we have discussed: effects of faulty mask alignment; differences between slit width and step size; excessive gaps between encoding and blocking masks, distortions produced by continuously moving, rather than stepped, masks, and so on.

Item 3, ii) is a complicated problem that turned out to have an analytical solution, described in a previously published paper, written before the grant award had actually been made.

"Errors in Hadamard Transform Spectroscopy or Imaging Caused by Imperfect Masks" M.H. Tai, M. Harwit, and N.J.A. Sloane, Appl.Opt. 14, 2678 (1975).

Errors listed under item 2 of section 1 also are treated in paper 4 (again, see Appendix). We were able to treat the effect of noise spikes, random noise and sinusoidally shaped noise froms. From these other types of noise can be analyzed by means of Fourier techniques. The question of obtaining more data points than are strictly necessary is one whose mathematical treatment bears the name of 'theory of singular weighing designs'. Again, the detailed analysis is difficult to summarize, but the mathemetical treatment is given in Appendix A, and we provide results there on optimum mask patterns that provide the best signal to noise ratios.

Finally, turning to item 1, the theory of singular designs also applies here, and again optimum mask patterns can be selected for given conditions. Improvements in signal to noise ratio can in fact be obtained by constraining some of the values of the spectrum's end elements to take on a constant known brightness. Here the precise value of the brightness enters not directly into the signal-to-noise ratio. It affects system performance more immediately in the requirements on dynamic range. By choosing the frame brightness to equal the mean brightness of the entire spectrum or scene, the dynamic range required by the data handling system can be made small, and this facilitates data handling.

3. Publications

1. Fourier and Hadamard transform spectrometers: a limited comparison, Ming Hing Tai and Martin Harwit, Applied Optics 15, 2664 (1976).
2. Fourier and Hadamard transform spectrometers: a limited comparison II, Martin Harwit and Ming Hing Tai, Applied Optics, 16, 3071 (1977).
3. Hadamard transform techniques in spectroscopy and imaging spectroscopy, M. Harwit, in "Transform Techniques in Chemistry" edited by Peter Griffiths, Plenum Press 1978.
4. Errors in Hadamard transform optics, N. J. A. Sloane, Martin Harwit and Ming Hing Tai, to be published, Applied Optics, 1978
5. "A practical Hadamard transform spectrometer for astronomical application," Ming Hing Tai, Cornell University Ph.D. thesis and Center for Radiophysics and Space Research Report, CRSR 655 (208 pages).

4. Ph.D. Thesis Work

A graduate student, Mr. Ming Hing Tai, was partially supported by the present contract while working on his Ph.D. thesis. His thesis work dealt with the performance of Hadamard transform spectrometers and the astronomical results that could be obtained with them, particularly in planetary observations. The thesis appeared as a Center for Radiophysics and Space Research (Cornell University) report CRSR 655, listed above.

Part of this 208 page compendium involves procedures for dealing with noise and systematic errors.

5. Personnel

Prof. M. Harwit, Principal Investigator

Mr. M. H. Tai, Graduate Research Assistant

Mr. G. Melnick, Graduate Reserach Assistant

Mr. F. Dain, Undergraduate Reserach Assistant

Mr. G. Stasavage, Research Technician

6. Budget Status

No uncommitted funds, as of November 20, 1977.

Appendix

Distortion in Hadamard Transform Optics

N. J. A. Sloane, Martin Harwit and Ming Hing-Tai

to be published in Applied Optics 1978

Distortion in Hadamard Transform Optics

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ABSTRACT

This paper analyzes the distortion in Hadamard transform optical instruments caused by moving masks, incorrect mask alignment, missing measurements, diffraction, etc., and describes techniques for reducing or eliminating this distortion. In a great many cases the behavior of the instrument can be characterized by a single matrix equation of the form $\eta = TWa + e$, where the components of η are the measurements, T is a matrix characterizing the instrument, W specifies the mask configurations, a is a vector containing the unknown spectral intensities, and the components of e are small random errors.

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Distortion in Hadamard Transform Optics

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1. Introduction

Hadamard transform optical systems have been widely studied during the past few years. By encoding the light with properly designed masks these instruments benefit from the so-called *multiplex advantage* and achieve an increased signal-to-noise ratio. A variety of multiplexing instruments have been described in the literature, designed to improve images, spectra or both.¹⁻¹⁰

Any optical technique that leads to improved performance is likely to have its limitations. Some of these limitations may be inherent in the design of the system, while others may be caused by practical difficulties with the apparatus. In either case once such a limitation is recognized ways of working with or around it can be examined. The purpose of this paper is to describe a number of sources of error encountered in Hadamard transform optical instruments, to give techniques for reducing or eliminating these errors where such techniques have already been worked out, and to suggest some possible ways of dealing with errors that have not yet been fully studied.

We shall mainly describe errors occurring in singly encoded spectrometers, although most of the discussion will apply equally well to imagers. Similar errors occur in doubly encoded systems but have not yet been extensively investigated.

2. Description of Instrument and Spectrum: Simplest Case

In this section we consider a multiplexing spectrometer in its simplest mode of operation. We give a mathematical description of the instrument under the assumption that no errors occur except random fluctuations in the detector readings. Subsequent sections will consider various departures from this ideal behavior.

2.1 Description of Instrument

Let us consider then a spectrometer with a single narrow entrance slit and with a mask in the exit focal plane.^{3,5} It is assumed that the entrance slit is sufficiently narrow so as to be always filled by the incoming radiation. The input to the spectrometer can therefore be completely specified by giving the frequency distribution of the radiation. It is most convenient to express this as a function $F(\nu)$, say, of the wavenumber ν . We further assume that over the operating range in which we are interested the optical system between the entrance slit and the exit mask behaves like a linear system.¹¹ In particular if the input consists of monochromatic

light at wavenumber ν_0 and with unit intensity, so that $F(\nu) = \delta(\nu - \nu_0)$, then the distribution of intensity along the exit focal plane can be written as

$$H(\nu - \nu_0).$$

(We assume that distances along the exit focal plane have been calibrated in terms of wavenumber.) Thus $H(\nu - \nu_0)$ is the instrument's spectral display for monochromatic input at wavenumber ν_0 . $H(\nu - \nu_0)$ is called the *impulse response* or *point spread function*^{1,2} of the instrument. Because of aberrations and diffraction $H(\nu - \nu_0)$ can never be concentrated at a single point, no matter how narrow the entrance slit is. By superposition the distribution of intensity along the exit focal plane produced by an arbitrary input $F(\nu)$ is given by

$$G(\nu) = \int_{-\infty}^{\infty} F(\nu_0) H(\nu - \nu_0) d\nu_0. \quad (1)$$

In other words $G(\nu)$ is the spectral display along the exit focal plane if the input spectrum is $F(\nu)$.

An ideal instrument, free of aberrations and diffraction, would have impulse response $H(\nu - \nu_0) = \delta(\nu - \nu_0)$ and (from (1))

$$G(\nu) = F(\nu).$$

Thus the linear spectral display $G(\nu)$ in the exit focal plane would be identical to the spectral distribution of the input. But in any real instrument $G(\nu) \neq F(\nu)$.

Of course strictly speaking we should write G as a function $G(x)$ of the position, x , along the exit focal plane, rather than as a function of wavenumber. But the latter notation is simpler and more suggestive.

Suppose a single exit slit is used, extending say from wavenumber ν_1 to wavenumber ν_2 along the exit focal plane. The detector output is given by the integral of the exit plane distribution over the slit area:

$$\begin{aligned} \eta(\nu_1, \nu_2) &= \int_{\nu_1}^{\nu_2} G(\nu) d\nu + e \\ &= \int_{-\infty}^{\infty} F(\nu_0) d\nu_0 \int_{\nu_1}^{\nu_2} H(\nu - \nu_0) d\nu + e \\ &= \int_{-\infty}^{\infty} F(\nu_0) S(\nu_0; \nu_1, \nu_2) d\nu_0 + e, \end{aligned} \quad (2)$$

where e is the error due to detector noise, and $S(\nu_0; \nu_1, \nu_2)$ is defined by

$$S(\nu_0; \nu_1, \nu_2) = \int_{\nu_1}^{\nu_2} H(\nu - \nu_0) d\nu \quad (3)$$

$S(\nu_0; \nu_1, \nu_2)$ is called the *slit function*^{1,2} of the instrument; it is the total intensity of radiation passing through a slit extending from ν_1 to ν_2 in the exit focal plane, when the input is monochromatic radiation at wavenumber ν_0 .

Figures 1(a) and 1(b) show two examples of $H(\nu - \nu_0)$. Figure 1(a) corresponds to an input slit which is wide enough for diffraction to be ignored, while Fig. 1(b) shows the opposite extreme when diffraction effects dominate and $H(\nu - \nu_0) = \text{sinc}^2(\nu - \nu_0)$, where $\text{sinc } x = (\sin x)/x$. The main lobe of $H(\nu - \nu_0)$ has width $2d_1$ in Fig. 1(a), or $2d_2$ in Fig. 1(b). The corresponding slit functions $S(\nu_0; \nu_1, \nu_2)$ are shown in Figs. 1(c), 1(d), assuming the exit slit has width $\nu_2 - \nu_1$ equal to $2d_1$ or $2d_2$ respectively. For example in Fig. 1(c) when $\nu_0 = (\nu_1 + \nu_2)/2$ the exit slit extends from $\nu_0 - d_1$ to $\nu_0 + d_1$ and gathers all the radiation, while

if $\nu_0 = (\nu_1 + \nu_2)/2 \pm d_1$ only half of the radiation passes through the slit.

Equation (2) is the fundamental equation that relates the output of the detector to the input spectrum $F(\nu)$ when the exit slit extends from $\nu = \nu_1$ to $\nu = \nu_2$.

2.2 Description of Spectrum

It is customary when measuring a spectrum with a dispersion instrument to divide the spectrum into n intervals, estimate the intensity in each interval, and then join the estimates by straight line segments. The result is a *piecewise linear curve* as shown in Fig. 2. Any spectrum can be approximated as closely as we please by such a curve, if n is taken large enough. We shall therefore *start off* with the assumption that the input spectrum $F(\nu)$ is a piecewise linear curve as shown in Fig. 2: the spectrum is divided into n segments $\theta'_0\theta'_1, \theta'_1\theta'_2, \dots, \theta'_{n-1}\theta'_n$, corresponding to positions $\theta_0\theta_1, \theta_1\theta_2, \dots, \theta_{n-1}\theta_n$ of the slits in the mask, and a_i is the intensity at the center of $\theta_i\theta_{i+1}$. (We assume that the scales have been chosen so that θ'_i and θ_i are related by an equation of the form $\theta'_i = \theta_i + A\theta_i^2 + \dots$, where A is a constant.)

The goal is to determine the n unknowns a_0, a_1, \dots, a_{n-1} as accurately as possible. In order to have an unambiguous expression for $F(\nu)$ in the first and last segments we assume that $F(\nu)$ is periodic outside the interval $\theta_0\theta_n$, as shown by the broken lines in Fig. 2. Thus we define $a_{-1} = a_{n-1}$ and $a_n = a_0$ as the intensities in the segments $\theta_{-1}\theta_0$ to the left of $\theta_0\theta_1$, and $\theta_n\theta_{n+1}$ to the right of $\theta_{n-1}\theta_n$, respectively. Analytically we are assuming that $F(\nu)$ is given by

$$F(\nu) = a_i + (a_{i+1} - a_i)(\nu - \theta_{i+1} + b)/2b \quad (4)$$

$$\text{for } \theta_{i+1} - b \leq \nu \leq \theta_{i+1} + b$$

where $2b = \theta_{i+1} - \theta_i$ is the length of each segment.

2.3 Operation as a Multiplexed Spectrometer

Let us consider what happens when the instrument is operated as a multiplexed spectrometer, by allowing light from several slits to fall simultaneously onto the detector during each measurement. Suppose there are n slit positions in all: $\theta_0\theta_1, \theta_1\theta_2, \dots, \theta_{n-1}\theta_n$. The pattern of open and closed slits is specified by means of the configuration matrix $W = (w_{ij})$, where $w_{ij} = 1$ if the j -th slit is open during the i -th measurement and $w_{ij} = 0$ otherwise, for $0 \leq i, j \leq n-1$. In statistical terminology W specifies a *weighing design*.¹³⁻¹⁵ The mechanical design of the instrument is simplified (as will appear in the next section) if W is either a (right) circulant matrix, with $w_{ij} = w_{i-j}$ or a left circulant matrix, with $w_{ij} = w_{j-i}$ (and subscripts taken modulo n in both cases). From now on we assume that W is a left circulant; similar results hold in the other case.

If the j -th slit is open during a particular measurement the total intensity of light passing through that slit is

$$\int_{\theta_j}^{\theta_{j+1}} G(\nu) d\nu \quad (5)$$

Since $G(\nu)$ is a linear function of $F(\nu)$ (Eq. (1)) and $F(\nu)$ is a linear function of the a_i (Eq. (4)) this intensity may be written as

$$\sum_{k=0}^{n-1} t_{j-k} a_k \quad (6)$$

where the t_{j-k} are constants characterizing the instrument. Thus t_{j-k} specifies the amount of radiation that is transferred through an open slit at position j , for unit intensity at the k -th spectral element. We assume that this is only a function of the separation $j-k$, independent of the actual values of j and k . To include more general situations (for example, misaligned masks - see Sections 3 and 4) we allow these constants to vary from measurement to measurement. Thus we assume that if the j -th slit is open during the i -th measurement, the total intensity of light passing through that slit (i.e. the contribution of this slit to the i -th detector

reading) is given by

$$\tau_j = \sum_{k=0}^{n-1} t_{j-k} a_k, \quad (7)$$

where again the t_{j-k} are constants depending on the particular instrument. The i -th detector reading η_i is given by the sum of all τ_j for which the corresponding slit is open:

$$\eta_i = \sum_{j=0}^{n-1} \tau_j w_{ij} + e_i \quad (8)$$

$$\begin{aligned} &= \sum_{k=0}^{n-1} \sum_{j=0}^{n-1} t_{j-k} w_{ij} a_k + e_i \\ &= \sum_{k=0}^{n-1} c_{ik} a_k + e_i \end{aligned} \quad (9)$$

where e_i is the error in the i -th reading and the constants c_{ik} are given by

$$\begin{aligned} c_{ik} &= \sum_{j=0}^{n-1} t_{j-k} w_{ij} \\ &= \sum_{j=0}^{n-1} t_{j-k} w_{r+j} \quad (\text{since } W \text{ is a left circulant}) \\ &= \sum_{r=0}^{n-1} t_{r-k} w_{r+k} \quad (\text{where } r = i+j-k) \\ &= \sum_{r=0}^{n-1} t_{r-k} w_{r,k} \end{aligned} \quad (10)$$

This is simpler in matrix notation. Let

$$\eta = (\eta_0, \eta_1, \dots, \eta_{n-1})^T,$$

$$a = (a_0, a_1, \dots, a_{n-1})^T,$$

$$e = (e_0, e_1, \dots, e_{n-1})^T$$

denote column vectors of measurements, unknowns, and errors, respectively. From (9) and (10) we obtain

$$\eta = Ca + e \quad (11)$$

$$= TWa + e, \quad (12)$$

where C is the matrix with (i,k) -th entry c_{ik} , and

$$T = \begin{bmatrix} t_0^0 & t_1^0 & t_2^0 & \dots & t_{n-1}^0 \\ t_0^1 & t_1^1 & t_2^1 & \dots & t_{n-2}^1 \\ \dots & \dots & \dots & \dots & \dots \\ t_0^{n-1} & t_1^{n-1} & t_2^{n-1} & \dots & t_0^{n-1} \end{bmatrix}. \quad (13)$$

T is called the *transfer matrix* of the instrument. In an ideal instrument in which distortion, aberrations and diffraction were negligible T would be the unit matrix I . In a real instrument, however, radiation that should be exiting through a given mask slit often spills over onto neighboring mask positions, so that diagonal elements of T are reduced in value and off-diagonal elements of the matrix grow. Since radiant energy incident on a detector cannot be negative, $0 \leq t_{j-k} \leq 1$. Conservation of radiant energy requires that the sum of the matrix elements in any row or column of T equals unity, unless dissipative processes such as absorption or

scattering play a role. If dissipation is significant these sums, $\sum_i t'_{i-k}$ or $\sum_k t'_{i-k}$, may be less than unity and greater than zero. A zero value for any one of these sums would imply a systematic blockage in the spectrometer -- either by design, or through faulty construction. If no light is lost then T satisfies

$$TJ = T^{-1}J = J$$

where J is a square matrix of ones.

Equations (11) and (12) are the basic equations describing the performance of this multiplexing spectrometer. They relate the unknowns a_0, \dots, a_{n-1} of Fig. 2 to the measurements $\eta_0, \dots, \eta_{n-1}$ via the matrix $C = (c_{ik})$ or via the matrices T and W . Here $W = (w_{ij}) = (w_{i,j})$ is the configuration matrix that describes which slits are open and which are closed during the measurements, and the transfer matrix T (Eq. (13)) characterizes the particular instrument but is independent of the multiplexing. The entries in T are determined by Eq. (7) and depend on the impulse response $H(\nu)$ and on the positions $\theta_0, \dots, \theta_{n-1}$. We shall see below that by suitably choosing the transfer matrix T in (12) we can use the same equation to describe the distortion introduced by a number of different sources of error.

2.4 Examples

(1): No Diffraction

As a first example we assume a wide entrance slit and no diffraction, so that the impulse response is as shown in Fig. 1(a). We take $d_1 = 1/2$, so that $H(\nu - \nu_0) = 1$ for $-1/2 \leq \nu - \nu_0 \leq 1/2$, $H(\nu - \nu_0) = 0$ elsewhere. We also assume that the exit slits have width 1, so that $\theta_{j+1} - \theta_j = 1$ and $b = 1/2$ in Eq. (4) and Fig. 2. The slit function $S(\nu_0; \theta_j, \theta_{j+1})$ is given in Fig. 1(c), or analytically by

$$S(\nu_0; \theta_j, \theta_{j+1}) = \begin{cases} \nu_0 - \theta_j + \frac{1}{2} & \text{if } \theta_j - \frac{1}{2} \leq \nu_0 \leq \theta_j + \frac{1}{2} \\ -\nu_0 + \theta_j + \frac{3}{2} & \text{if } \theta_j + \frac{1}{2} \leq \nu_0 \leq \theta_j + \frac{3}{2} \\ 0 & \text{otherwise} \end{cases} \quad (14)$$

If the j -th slit is open during the i -th measurement, the total intensity of light passing through that slit is found by substituting $S(\nu_0; \theta_j, \theta_{j+1})$ and $F(\nu)$ into Eq. (2), to give

$$\begin{aligned} \tau'_j &= \int_{\theta_j - 1/2}^{\theta_j + 1/2} (\nu_0 - \theta_j + \frac{1}{2}) (a_{j-1} + (a_j - a_{j-1})(\nu_0 - \theta_j + \frac{1}{2})) d\nu_0 \\ &\quad + \int_{\theta_j + 1/2}^{\theta_j + 3/2} (-\nu_0 + \theta_j + \frac{3}{2}) (a_j + (a_{j+1} - a_j)(\nu_0 - \theta_j - \frac{1}{2})) d\nu_0 \\ &= \frac{1}{6} (a_{j-1} + 4a_j + a_{j+1}) \end{aligned} \quad (15)$$

Therefore the constants t'_{j-k} in Eq. (7) are given by

$$\begin{aligned} t'_0 &= \frac{4}{6}, \quad t'_{-1} = t'_1 = \frac{1}{6}, \\ t'_j &= 0 \quad \text{otherwise} . \end{aligned}$$

(Note that in this instrument the t'_j are independent of i .) The transfer matrix T for this instrument is therefore

$$T = \frac{1}{6} \begin{bmatrix} 4 & 1 & 0 & \cdots & 0 & 1 \\ 1 & 4 & 1 & \cdots & 0 & 0 \\ 0 & 1 & 4 & \cdots & 0 & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\ 0 & 0 & 0 & \cdots & 4 & 1 \\ 1 & 0 & 0 & \cdots & 1 & 4 \end{bmatrix} . \quad (16)$$

Example (2): Full Diffraction

The second example assumes a very narrow entrance slit, so that the impulse response is

$$H(\nu - \nu_0) = 2 \operatorname{sinc}^2(2\pi(\nu - \nu_0)) , \quad (17)$$

as in Fig. 1(b) with $d_2 = 1/2$. (The constant 2 is chosen so that the total area under $H(\nu - \nu_0)$ is one). Again we assume the slits have width 1, $\theta_{r+1} - \theta_r = 1$, $b = 1/2$. Now the contribution of the j -th slit to the i -th measurement is (from Eqs. (2), (4), (17)):

$$\tau_j' = \sum_{r=0}^{n-1} \int_{\theta_{r+1}-1/2}^{\theta_{r+1}+1/2} \{a_r + (a_{r+1} - a_r)(\nu_0 - \theta_{r+1} + \frac{1}{2})\} d\nu_0 \cdot \int_{\theta_j}^{\theta_j+1} 2 \operatorname{sinc}^2 2\pi(\nu - \nu_0) d\nu$$

After some algebra this becomes

$$\tau_j' = t_0' a_j + (t_1' a_{j-1} + t_{-1}' a_{j+1}) + \dots$$

(i.e. Eq. (7)), where the coefficients t_r' are given by

$$t_r' = 2 \int_{-3/2}^{3/2} f(x) \operatorname{sinc}^2 2\pi(x - r) dx , \quad (18)$$

and

$$\begin{aligned} f(x) &= \frac{1}{2}(x + \frac{3}{2})^2 \text{ if } -\frac{3}{2} \leq x \leq -\frac{1}{2} , \\ f(x) &= \frac{3}{4} - x^2 \text{ if } -\frac{1}{2} \leq x \leq \frac{1}{2} , \\ f(x) &= \frac{1}{2}(x - \frac{3}{2})^2 \text{ if } \frac{1}{2} \leq x \leq \frac{3}{2} . \end{aligned} \quad (19)$$

Numerical evaluation of the integrals (18) shows that the transfer matrix T of this instrument is a symmetric circulant matrix with first row equal to

$$\begin{aligned} & (0.6667, 0.1482, 0.0080, 0.0031, 0.0017, \\ & 0.0010, 0.0007, 0.0005, 0.0004, 0.0003, \\ & 0.0003, 0.0002, \dots, 0.0017, 0.0031, \\ & 0.0080, 0.1482) \end{aligned}$$

(20)

correct to four decimal places. Notice the very slow decay of elements off the main diagonal, illustrating the pronounced spreading effect of this impulse response.

In general whenever $H(\nu)$ is an even function, i.e. satisfies $H(-\nu) = H(\nu)$, then T is a symmetric circulant — a matrix of the form

$$\begin{pmatrix} u_0 & u_1 & u_2 & \dots & u_3 & u_2 & u_1 \\ u_1 & u_0 & u_1 & \dots & u_4 & u_3 & u_2 \\ u_2 & u_1 & u_0 & \dots & u_5 & u_4 & u_3 \\ \vdots & \vdots & \vdots & \ddots & \vdots & \vdots & \vdots \\ u_1 & u_2 & u_3 & \dots & u_2 & u_1 & u_0 \end{pmatrix} . \quad (21)$$

2.6 Operation as a Monochromator

If the same instrument is operated as a conventional monochromator, n measurements are made, where in the i -th measurement a single exit slit extends from θ_i to θ_{i+1} . The basic equation is now

$$\eta = T\alpha + e, \quad (22)$$

obtained by setting $W = I$ in (12).

2.7 Recovery of Spectrum

To estimate the spectrum from the measurements we proceed as follows. We assume (again this is the ideal case) that e_i , the detector error in the i -th measurement, is a random variable with mean zero:

$$E\{e_i\} = 0, \quad (23)$$

that the errors in distinct measurements are uncorrelated:

$$E\{e_i e_j\} = 0, \quad i \neq j, \quad (24)$$

and that

$$E\{e_i^2\} = \sigma^2 \quad (25)$$

where the variance σ^2 is independent of the amount of radiation falling on the detector. Under these assumptions if (12) holds then the best estimate of the α 's is given by

$$\hat{\alpha} = W^{-1} T^{-1} \eta \quad (26)$$

$$= \alpha + W^{-1} T^{-1} e. \quad (27)$$

This is best in the sense of being that linear unbiased estimate which minimizes the average mean square error

$$\begin{aligned} \epsilon &= \frac{1}{n} \sum_{i=0}^{n-1} E\{(\hat{\alpha}_i - \alpha_i)^2\} \\ &= \frac{\sigma^2}{n} \text{Trace}\{(TW)^{-1}((TW)^{-1})^T\} \\ &= \frac{\sigma^2}{n} \times \{\text{sum of squares of elements of } (TW)^{-1}\} \end{aligned} \quad (28)$$

-- see Refs. 3,18. Then ϵ is a measure of the accuracy of the experiment, and should be made as small as possible.

2.8 Computing the Spectrum. Inverse Matrices

The spectrum is estimated by (26). Usually W^{-1} is known explicitly,³ but problems may arise in finding T^{-1} , since in general the inverse of say a 255×255 matrix is both difficult to compute and awkward to store.

If T is a circulant matrix, as is often the case, then the inverse is fairly easy to find and to store -- see below. If T is Toeplitz (constant along diagonals) but not a circulant, then the inverse can still be found but with more difficulty.¹⁹⁻²¹ But usually if T is not a circulant then the best method of computing the spectrum is not to attempt to find T^{-1} but to directly solve the system of equations

$$\eta = TW\alpha \quad (29)$$

for α . This is particularly straightforward in the important case when W^{-1} is known and T is a *band matrix*, i.e. the only nonzero entries in T are those within a fixed distance of the main

diagonal.²²

The following are two methods for finding the inverse of an invertible circulant matrix C with first row c_0, c_1, \dots, c_{n-1} . Method (I)^{16,17} often gives an explicit formula for the inverse: Let $c(x) = c_0 + c_1x + \dots + c_{n-1}x^{n-1}$. Find the unique polynomial $d(x) = d_0 + d_1x + \dots + d_{n-1}x^{n-1}$ such that $c(x)d(x) \equiv 1 \pmod{x^n - 1}$. Then C^{-1} is the circulant with first row d_0, \dots, d_{n-1} . Method (II) (which may be new) is an efficient algorithm for finding the inverse of a symmetric circulant with a computer: For large n the m -th entry d_m of the first row of the inverse approaches

$$\frac{1}{\pi} \int_0^\pi \frac{\cos m\theta d\theta}{c_0 + 2 \sum_{r=1}^{\infty} c_r \cos r\theta}. \quad (30)$$

In calculating (26) it is useful to remember that any two right circulants commute, i.e. satisfy $AB = BA$, any two left circulants commute, and any left circulant commutes with any symmetric right circulant.

Example (1) (cont.)

The inverse of (16) was found by Method (I). First define the sequence of integers b_0, b_1, b_2, \dots by

$$\begin{aligned} b_0 &= 1, & b_1 &= -4 \\ b_n &= -4b_{n-1} - b_{n-2}, & n &\geq 2. \end{aligned} \quad (31)$$

This is Sequence 1420 in Ref. 23. The first few terms are

$$\begin{array}{cccccc} b_0 & b_1 & b_2 & b_3 & b_4 & b_5 \\ 1 & -4 & 15 & -56 & 209 & -780 \end{array}$$

An explicit formula is

$$b_n = \frac{1}{6} \{ (3 - 2\sqrt{3})(-2 + \sqrt{3})^n + (3 + 2\sqrt{3})(-2 - \sqrt{3})^n \} \quad (32)$$

and also

$$b_{n-1}^2 - b_n b_{n-2} = 1, \quad n \geq 1. \quad (33)$$

Then the inverse of (16) is given by the symmetric circulant (21) with

$$u_0 = \frac{3b_{n-1}}{2b_{n-1} + b_{n-2} + 1}, \quad (34)$$

$$u_i = \frac{3(b_{n-i-1} + b_{i-1})}{2b_{n-1} + b_{n-2} + 1}, \quad i \geq 1. \quad (35)$$

(This may be verified by using Method (I); we omit the rather lengthy details.) For example when $n = 4$

$$T = \frac{1}{6} \begin{bmatrix} 4 & 1 & 0 & 1 \\ 1 & 4 & 1 & 0 \\ 0 & 1 & 4 & 1 \\ 1 & 0 & 1 & 4 \end{bmatrix}, \quad T^{-1} = \frac{1}{4} \begin{bmatrix} 7 & -2 & 1 & -2 \\ -2 & 7 & -2 & 1 \\ 1 & -2 & 7 & -2 \\ -2 & 1 & -2 & 7 \end{bmatrix}.$$

For large n the second term in (32) dominates and we can approximate u_i by

$$u_i \approx u'_i = (-1)^i \sqrt{3}(2 - \sqrt{3})^i \quad (36)$$

for $0 \leq i \leq n/2$, independent of n . The approximation is very accurate, in fact

$$|u_i - u'_i| < 10^{-5} \quad \text{if } n \geq 20.$$

Table I gives the values of u'_i , correct to three decimal places. Note that u'_i approaches 0 rapidly as i increases. Also

$$\det T \approx \left(\frac{2 + \sqrt{3}}{6} \right)^n. \quad (37)$$

Table I
First Row of Inverse Matrix of (16)

i	u'_i
0	1.732
1	-0.464
2	0.124
3	-0.033
4	0.009
5	-0.002
6	0.001

Therefore for example if the instrument is being operated as a monochromator, so that Eq. (22) applies, the best estimate of the a_i 's is given by

$$\hat{a} = T^{-1}\eta. \quad (38)$$

and for $n \geq 20$ we can write Eq. (38) as

$$\hat{a}_i = u'_0\eta_i + u'_1(\eta_{i-1} + \eta_{i+1}) + \cdots + u'_6(\eta_{i-6} + \eta_{i+6})$$

with an error of less than 10^{-3} , where the u'_i are given by Table I, and subscripts are read modulo n . The average mean square error ϵ for this mode of operation can be found from Eqs. (28), (36) and is

$$\epsilon \approx 2\sqrt{3} \sigma^2. \quad (39)$$

Incidentally, if the correction matrix T^{-1} is not applied, the result is a coarsening of the spectrum. Consider an input spectrum with a single sharp line, corresponding to $a = (0, 0, 0, 1, 0, 0, \dots)$ as shown in Fig. 3(a). Then from Eq. (22), ignoring e ,

$$\eta = Ta = (0, 0, 1/6, 2/3, 1/6, 0, 0, \dots).$$

If we ignore T^{-1} and estimate \hat{a} by η (instead of (38)) the resulting spectrum is shown in Fig. 3(b).

Suppose the same instrument is operated as a multiplexed spectrometer, so that Eq. (12) applies. The argument used in the Appendix of Ref. 15 can be modified to show that for any (0,1)-matrix W we have

$$\epsilon \geq \frac{10.34 \sigma^2}{n}, \quad n \text{ large}, \quad (40)$$

(using Eq. (37)). If W is an S-matrix^{3, 15} then it can be shown that

$$\epsilon \approx \frac{13.86 \sigma^2}{n}, \quad n \text{ large}, \quad (41)$$

which is reasonably close to (40) and may well be the smallest ϵ that can be achieved. Comparing (39) with (41) we see that

$$\frac{\text{m.s.e. with multiplexing}}{\text{m.s.e. without multiplexing}} \approx \frac{4.00}{n} \quad (42)$$

and in fact for this choice of W it can be shown that (42) holds for any matrix T .

Example (2) (cont)

The inverse of (20) was found by Method (II), and is a circulant with first row which approaches

$$(1.6683, -0.3820, 0.0705, -0.0183, 0.0017, \\ -0.0017, -0.0006, -0.0006, -0.0004, -0.0003, \dots) \quad (43)$$

correct to four decimal places, for large n . Again the elements off the main diagonal decay very slowly. Without multiplexing the average mean square error for this instrument is

$$\epsilon \approx 2.93 \sigma^2,$$

while multiplexing with an S-matrix gives

$$\epsilon \approx \frac{11.74 \sigma^2}{n},$$

and again (42) holds.

Of course for any instrument we can ask the following question. The impulse response $H(\nu)$ determines the transfer matrix T , as in Eq. (13). Then which configuration matrix W minimizes the average mean square error (28)? In Ref. 15 we studied the performance of various W 's in the very special case when T is the identity matrix (thus ignoring the spreading effect of $H(\nu)$). The more general question for an arbitrary $H(\nu)$ or T remains unsolved.

3. Errors Occurring with a Mask which is Moved in Steps

This section deals with some classes of errors that can occur with a mask which is stepped to the next position between measurements. Most of these can be described by a suitable modification of the transfer matrix T in the basic equation (12). These errors may also arise in instruments with a continuously moving mask (see § 4). The resulting T is then the product of the matrices described in §4 and those below. The echo effects produced by slits which are uniformly too narrow (or too wide) were described by Tai et al.²⁴

3.1 Faulty Mask Alignment

Sometimes the position of an encoding mask is systematically displaced by some small distance δ from its correct position (see Fig. 4). Plankey et al.⁹ have described the effects produced by deliberate misalignment of a spectral encoding mask by 0.5 and 2.5 slit widths. As one would expect, the main result is an apparent spectral shift by 0.5 and 2.5 resolution elements, respectively, and an increase in spectral noise. The exact form of this noise may be determined as follows. We assume that δ is small compared to the slit width. Then the contribution to the i -th detector reading from light passing through the j -th slit is now given by

$$\tau_j' = \int_{\lambda_j - \delta}^{\lambda_{j+1} - \delta} G(\nu) d\nu \quad (44)$$

instead of Eq. (5). For the instrument of Example (1), this becomes

$$\begin{aligned}
\tau_j' = & \int_{\nu_j - 1/2 - \delta}^{\nu_j - 1/2} (-\theta_j + \frac{1}{2} + \delta + \nu_0) \{a_{j-2} + (a_{j-1} - a_{j-2})(\nu_0 - \theta_j + \frac{3}{2})\} d\nu_0 \\
& + \int_{\nu_j - 1/2}^{\nu_j + 1/2 - \delta} (-\theta_j + \frac{1}{2} + \delta + \nu_0) \{a_{j-1} + (a_j - a_{j-1})(\nu_0 - \theta_j + \frac{1}{2})\} d\nu_0 \\
& + \int_{\nu_j + 1/2 - \delta}^{\nu_j + 1/2} (\theta_j + \frac{3}{2} - \delta - \nu_0) \{a_{j-1} + (a_j - a_{j-1})(\nu_0 - \theta_j + \frac{1}{2})\} d\nu_0 \\
& + \int_{\nu_j + 1/2}^{\nu_j + 3/2 - \delta} (\theta_j + \frac{3}{2} - \delta - \nu_0) \{a_j + (a_{j+1} - a_j)(\nu_0 - \theta_j - \frac{1}{2})\} d\nu_0 \\
= & \frac{1}{6} (a_{j-1} + 4a_j + a_{j+1}) + \frac{1}{2} \delta (a_{j-1} - a_{j+1}) \\
& + \frac{1}{2} \delta^2 (a_{j-1} - 2a_j + a_{j+1}) + \frac{1}{6} \delta^3 (a_{j-2} - 3a_{j-1} + 3a_j - a_{j+1})
\end{aligned} \tag{45}$$

Hence the transfer matrix T is an asymmetric circulant with first row equal to

$$\frac{1}{6} (4 - 6\delta^2 + 3\delta^3, (1 - \delta)^3, 0, \dots, 0, \delta^3, 1 + 3\delta + 3\delta^2 - 3\delta^3). \tag{46}$$

Thus a sharp spectral line is dispersed asymmetrically into the adjoining resolution elements. Furthermore a sharp line at one end of the spectrum affects elements at the opposite end of the spectrum.

To deal with this problem the mask should first be repositioned so as to minimize the misalignment. The remaining error can then be removed by using the correct T^{-1} in Eq. (26).

3.2 Differences Between Slit Width and Step Size

When a mask contains large number of slits, a similarly large number of steps must be taken in order to move the encoding mask from one extreme position to the other. Over such a range of steps, small differences between the step size and the encoding slit width can accumulate. For example, a typical spectral mask might have a slit width of 0.1 mm and be 255 slits wide. Over this total width of 2.55 cm, the mask must be stepped with sufficient precision that its final position is precisely one step short of cycling into the initial configuration. Thus the systematic error in the mask motion has to be less than 4×10^{-3} mm per step, for the final position to be within one-tenth of a slit width from the intended location.

If such precision cannot be attained then sharp spectral lines will take on a broadened appearance. Figure 5 shows a computer simulation of this effect. The analysis of this distortion is very similar to that in Sec. 3.1.

Suppose the initial position of the mask is correct, and the final position is $n\Delta$ behind the correct position. Then Eq. (5) must be replaced by

$$\tau_j' = \int_{\nu_j - i\Delta}^{\nu_j + 1 - i\Delta} G(\nu) d\nu. \tag{47}$$

The transfer matrix T for the instrument of Example (1) has i -th row equal to

$$\begin{aligned}
& (0, \dots, 0, (i\Delta)^3, 1 + 3i\Delta + 3(i\Delta)^2 - 3(i\Delta)^3, \\
& 4 - 6(i\Delta)^2 + 3(i\Delta)^3, (1 - i\Delta)^3, 0, \dots, 0)
\end{aligned} \tag{48}$$

with the entry $4 - 6(i\Delta)^2 + 3(i\Delta)^3$ on the main diagonal.

3.3 Excessive Gap Between Encoding and Blocking Masks

If the blocking mask and encoding mask are mounted too close together, there is a danger that one of the opaque elements of the encoding mask may catch the edge of the blocking mask. To avoid this there must be a gap between the two masks. If this gap is too large, however, some of the radiation that strikes the blocking mask obliquely will pass into an encoding slit that actually was meant to be obscured by the blocking mask and therefore inaccessible to any incident radiation. Then a spectral line imaged on either extreme end of the blocking mask will produce a false spectral line at the opposite end of the spectrum. This effect is most pronounced in fast optical systems (with small focal ratios).

As long as the encoding (rather than the blocking) mask is placed in the plane where the sharpest focus is obtained, this effect is not very serious because it only affects the extreme elements at the two ends of the spectrum. In contrast the mask misalignment mentioned in §3.2 produces distortion over the entire spectrum.

3.4 Nonlinearities

Any of the versions of Eq. (12) derived in this paper will only be valid as long as there is a linear relation between the input and the detector readings. Eventually this linear relation will fail, as the number of spectral elements increases, or as the resolution is increased, for example. Equation (12) must then be replaced by a system of simultaneous nonlinear equations. Although numerical techniques for solving systems of nonlinear equations are available,^{25, 26} little work has been done so far to determine the exact form of the equations that will be needed to replace (12).

Some nonlinearities can be compensated for very simply:

(a) A nonlinear response of the detector to increasing amounts of radiation can be eliminated through judicious calibration which permits nonlinear (compensating) scaling of the actual detector output before further data processing is undertaken.

(b) A nonlinear wavelength or wavenumber response -- when $G(\nu)$, the spectral display function along the exit plane, is not a linear function of ν -- can be compensated for by suitably plotting the final spectrum in a way that takes the actual wavelength calibration of the instrument into account.

(c) However, if the detector response at one wavenumber is a nonlinear function of the other wavenumbers present, compensation becomes very difficult, and the best procedure may be to solve for the spectrum iteratively, by successive approximations.

4. Distortion Introduced by a Continuously Moving Mask

In this section we analyze the distortion introduced when the slits move continuously across the exit focal plane instead of being discretely stepped between measurements. (An alternative analysis, which applies also to imagers, has been given by Gottlieb.¹ However, the matrix approach given here seems simpler.) We shall see that in many cases the basic equation (12) still holds, with T given by (13), where the entries t'_{j-k} are obtained from Eq. (7). Again the spectrum is estimated by $\hat{a} = W^{-1}T^{-1}\eta$, and the effect of T^{-1} is (theoretically at least) to eliminate the distortion caused by the moving slits.

4.1 Derivation of Basic Equation

As mentioned in §2.3 we assume that the configuration matrix is a left circulant:

$$W = \begin{bmatrix} w_0 & w_1 & \dots & w_{n-1} \\ w_1 & w_2 & \dots & w_0 \\ w_2 & w_3 & \dots & w_1 \\ \vdots & \vdots & \ddots & \vdots \\ w_{n-1} & w_0 & \dots & w_{n-2} \end{bmatrix}. \quad (49)$$

In many cases the S-matrices^{3, 15} can be arranged so as to have this form. An example with

$n = 7$ is given in Eq. (50).

$$s_7 = \begin{bmatrix} 1 & 1 & 1 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 & 1 \\ 1 & 0 & 1 & 0 & 0 & 1 & 1 \\ 0 & 1 & 0 & 0 & 1 & 1 & 1 \\ 1 & 0 & 0 & 1 & 1 & 1 & 0 \\ 0 & 0 & 1 & 1 & 1 & 0 & 1 \\ 0 & 1 & 1 & 1 & 0 & 1 & 0 \end{bmatrix}$$

(50)

Then instead of using n separate masks of length n -- one for each row of (49) -- we can use one long mask of length $2n$ which is moved continuously across the exit plane. Figure 6 shows the mask of length 14 corresponding to (50). The extra half-slit at each end of the mask is necessary to avoid errors at the ends of the spectrum. If $w_{n-1} = 0$ in (49) the extra half-slits are opaque. The mask is periodic with period n .

Ideally the slits have width $2b = \theta_{j+1} - \theta_j$. Let T_0 be the time for each measurement. When operated correctly the mask should move at a constant velocity $2b/T_0$, so as to move one slit width in the time taken to make a measurement.

In general, as in the previous section we let τ_j^i denote the contribution towards the i -th detector reading from light passing through an open slit which crosses the j -th segment θ_j, θ_{j+1} (see Fig. 7). We may write

$$\tau_j^i = \int_{\nu_{i,j}}^{\nu_{i+1,j+1}} l(\nu) G(\nu) d\nu, \quad (51)$$

where $\nu_{i,j}, \nu_{i,j+1}$ are the endpoints of this slit at the start of the i -th measurement, $\nu_{i+1,j+1}, \nu_{i+1,j+2}$ are the endpoints of the slit at the conclusion of the i -th measurement, and $l(\nu)$ is the length of time during which light at wavenumber ν can pass through the slit. Thus $l(\nu)$ can be read off Fig. 7: it is the length of the intersection of a vertical line at position ν with the shaded region which is the path swept out by the slit. By combining (8) and (51) we again get the basic equation (12).

4.2 An Ideal Case

Figure 8 shows the ideal case, when the mask has the correct dimensions, is properly aligned, and moves at the right speed. In this case

$$I(\nu) = (\nu - \theta_1 + b) T_0 / 2b, \quad \theta_1 - b \leq \nu \leq \theta_1 + b.$$

$$I(\nu) = (\theta_{n+1} + b - \nu) T_0 / 2b, \quad \theta_{n+1} - b \leq \nu \leq \theta_{n+1} + b. \quad (52)$$

As an example, suppose the instrument of Example (1) is being operated in this ideal mode. It follows from Eqs. (4), (51), (52) that

$$\tau'_j = \frac{1}{384} (a_{j-2} + 76a_{j-1} + 230a_j + 76a_{j+1} + a_{j+2}) \quad (53)$$

Since this is of the same form as Eq. (7) we see that (12) still holds, where the transfer matrix T is now a circulant with first row equal to

$$\frac{1}{384} (230, 76, 1, 0, 0, \dots, 0, 0, 1, 76). \quad (54)$$

Actually this is an approximation. Since the input spectrum $F(\nu)$ is in fact not periodic as we have assumed, Eq. (53) must be modified for $i = 0, 1, 2, n-3, n-2, n-1$, and the top right-hand and bottom left-hand corners of T are missing. Thus T is a Toeplitz¹⁹ matrix, not a circulant. However, it seems worthwhile changing T to a circulant and accepting the resulting distortion in the ends of the spectrum, in order to obtain a matrix which has a manageable inverse. The inverse of the circulant (54) is easily obtained, e.g. by Method (II) of §2.8. The inverse is a circulant with first row which approaches

$$(2.213, -0.826, 0.299, -0.108, 0.039, \\ -0.014, 0.006, -0.002, 0.001, 0, 0, 0, \dots) \quad (55)$$

correct to three decimal places, as n increases. For $n = 30$ (55) is already valid to this order of accuracy.

Again the effect of ignoring T is to broaden the spectrum, but now the broadening is more pronounced than in Fig. 3.

4.3 If the Mask Velocity or Slit Width is Wrong

It is not difficult to modify the above analysis to determine the distortion introduced if the mask is misaligned or if the mask velocity or slit widths are incorrect. It can be shown that, as long as each slit has the same (possibly incorrect) width, Eqs. (12) and (13) still hold. E.g. for the instrument of Example (1) operating with a misaligned mask moving with incorrect velocity we find that

$$\tau'_j = t'_2 a_{j-2} + t'_1 a_{j-1} + t'_0 a_j + t'_{-1} a_{j+1} + t'_{-2} a_{j+2} \quad (56)$$

for certain constants t'_i (again with suitable modifications for the first and last few measurements).

It does not seem worthwhile giving further examples, since each instrument will have its own transfer matrix T to be used in Eq. (12), and T should be determined when the instrument is calibrated.

4.4 Rotating Two-Dimensional Masks

The above effects also occur in imagers,^{6,10} when the encoding mask may be a two-dimensional mask mounted on a rotating disk. For example, a slight radial eccentricity will cause an effective radial motion of the mask. The result of this run-out should be a superposition of three effects. There should be a widening of the image along the radial direction, similar to the effect discussed in §3.2 above. The overall appearance of the image should become noisier, because of the continuous motion, as in §§4.1, 4.2. Thirdly, there should be a single cycle of roughly sinusoidal change in intensity along the unfolded two-dimensional chain of elements, somewhat similar to an effect described in §5 below. This is because the eccentric motion of the mask will move an opaque portion of the wheel into the open frame in the blocking mask, reducing the intensity during that portion of the cycle.

4.5 Imaging a Moving Source

If one of the sources being observed moves slightly during the measurements, the results will be similar to those described in §3.2 and §4.4, depending on whether the motion is along the direction of the mask's motion or across it. By and large, the result of a small displacement, of the order say of one spatial element during the frame time (or complete cycling time) of the mask will be to stretch the image along the direction of motion, and to add noise to the entire scene.

5. Effect of Drift in Background Level

If the intensity of background radiation incident on a spectrometer or imaging system varies during a spectral run, the derived spectrum will necessarily be affected. Suppose for example that the infrared spectrum of a star is to be observed. If the foreground atmospheric emission drifts during the course of the spectral run the instrument will record a corresponding drift, and the final spectrum obtained for the star will contain a component which can be directly attributed to the drift. Because the spectrum is estimated by a linear operation (Eq. (26)), the drift component is simply added to the true stellar spectrum. In this section we analyze the magnitude of this effect for several different types of drift, assuming for simplicity that $T = I$. (The analysis in the general case is much more complicated and we do not go into it here.)

Let η denote the vector of measurements with no drift present (Eq. (12)), let d_j be the drift in the j -th measurement,

$$d = (d_0, d_1, \dots, d_{n-1})^T,$$

and let

$$\eta' = \eta + d \quad (57)$$

be the actual measurements. The estimate of the spectrum is

$$\begin{aligned} \hat{a} &= W^{-1} \eta' \\ &= W^{-1} \eta + W^{-1} d, \end{aligned} \quad (58)$$

and we wish to analyze the drift component

$$D = W^{-1} d. \quad (59)$$

For concreteness we take W to be a symmetric, left circulant S -matrix,^{1,15} with

$$W^{-1} = S^{-1} = \frac{2}{n+1} (2S - I). \quad (60)$$

Case (i) = Constant Offset

Suppose

$$d = \delta 1 \quad (61)$$

where $1 = (1, 1, \dots, 1)^T$. Then

$$D = \frac{2\delta}{n+1} 1. \quad (62)$$

Case (ii): Sinusoidal Drift

Suppose

$$d_j = \cos 2\pi a(j - \phi), \quad j = 0, 1, \dots, n-1, \quad (63)$$

where a, ϕ are constants, ϕ being a phase shift. The k -th component of the drift is

$$D_k = \sum_{j=0}^{n-1} \xi_j \cos 2\pi a(j - k - \phi) \quad (64)$$

where

$$(\xi_0, \xi_1, \dots, \xi_{n-1}) \quad (65)$$

is the first row of W^{-1} . Simplifying (64) we obtain

$$D_k = A \cos 2\pi a(k + \phi - \beta), \quad (66)$$

where the amplitude is given by

$$\begin{aligned} A^2 &= \frac{4n}{(n+1)^2} + 2 \sum_{i < j} \xi_i \xi_j \cos 2\pi a(i-j) \\ &\approx \frac{4}{n} \end{aligned} \quad (67)$$

for large n , since (65) is a pseudo-random sequence,²⁷ and the phase shift β is given by

$$\tan(2\pi a\beta) = \frac{\sum_{j=0}^{n-1} \xi_j \sin 2\pi a_j}{\sum_{j=0}^{n-1} \xi_j \cos 2\pi a_j}. \quad (68)$$

We conclude that a sinusoidal drift voltage adds to the true spectrum a sinusoid of the same frequency as the drift, with amplitude multiplied by $2/\sqrt{n}$, and with a phase shift given by (68). Notice that a sinusoidal drift affects the spectrum more strongly than a constant offset of equal amplitude.

Case (iii): A Single Noise Spike

Suppose $d_k = \delta$, $d_j = 0$ for $j \neq k$. Then $\{D_j\} = \{\delta \xi_{j+k}\}$ is a pseudo-random sequence of r.m.s. value

$$\left(\frac{1}{n} \sum_j D_j^2 \right)^{1/2} \approx 2\delta/n \quad (69)$$

Thus a single noise spike of amplitude δ produces fairly random noise over the whole spectrum with r.m.s. value $2\delta/n$.

Case (iv)

Conversely, if $d_j = \delta \xi_{j+k}$, D has a single spike of amplitude δ in the k -th component.

These results have been collected in Table II, where the different drifts have been normalized so as to have unit amplitude. The last line of the table shows the improvement in signal-to-noise ratio obtained by considering random noise, obtained from Eq. (42). It is interesting that this is approximately equal to the improvement obtained when the noise is coherent. Only a constant offset in the data gives rise to a considerably smaller mean square error in the final spectrum.

Similar results are obtained in Fourier spectroscopy. A noise spike in the data will again produce a wide distribution of error signals in the final spectrum. A sinusoidal noise component, on the other hand, will produce a single spectral line at a frequency corresponding to that of the noise. (The analogous noise in the Hadamard instrument is that considered in Case (iv).)

6. Singular Designs

Table II
Noise in Spectrum Produced by Drift of
Unit Amplitude in the Data

Type of drift	R.m.s. drift in data	Noise in spectrum	R.m.s. spectrum noise	Ratio of r.m.s. spectrum noise to r.m.s. drift
Constant offset	1	Constant offset	$2/n$	$2/n$
Sinusoid	$1/\sqrt{2}$	Sinusoid	$\sqrt{2}/n$	$2/\sqrt{n}$
Single spike	$1/\sqrt{n}$	Eq. (69)	$2/n$	$2/\sqrt{n}$
Random	-	Random	-	$2/\sqrt{n}$

6.1 More Measurements Than Unknowns

It is sometimes desirable to design the experiment so that the number of measurements, p , exceeds the number of unknowns, n . The purpose in doing this is to ensure that in case some measurements are lost (for example if a cloud passes over an observatory during an astronomical observation) the spectrum can still be calculated.²⁸

The basic equation (12) relating the vector of measurements $\eta = (\eta_0, \eta_1, \dots, \eta_{p-1})^T$ to the vector of unknowns a becomes

$$\eta = TWa + e, \quad (70)$$

where W is a $p \times n$ (0,1) configuration matrix describing the experiment, and T is a $p \times p$ transfer matrix characterizing the instrument. The (i,j) -th entry of T is t_{ij} , $0 \leq i, j \leq p-1$ (see Eq. (7)).

The best estimate for a is now^{3,18}

$$\hat{a} = (TW)^+ \eta, \quad (71)$$

where $(TW)^+$ is the generalized inverse, given in this case by

$$(TW)^+ = (W^T T^T TW)^{-1} (TW)^T. \quad (72)$$

The average mean square error is then³

$$\epsilon = \frac{\sigma^2}{n} \text{Trace}\{(W^T T^T TW)^{-1}\}. \quad (73)$$

A mask of this type suitable for use in a spectrometer can be obtained by taking W to be the first n columns of a $p \times p$ circulant S -matrix. This can be accomplished by means of a blocking mask which only exposes a frame of n slits at a time.

6.2 More Unknowns Than Measurements

Suppose an experiment has been designed to make n measurements in order to determine a spectrum with n unknown components, via Eqs. (12), (26), but is prematurely terminated after only $p < n$ measurements have been made. In some cases it is still possible to say something about the unknown spectrum. For example, suppose $W = S_n$ and $T = I$. Then $\eta = S_n a$, and the sum of all n measurements is equal to $(n+1)/2$ times the sum of the unknown spectral components:

$$\sum_{j=1}^n \eta_j = \frac{1}{2}(n+1) \sum_{i=1}^n a_i. \quad (74)$$

Suppose the last measurement η_n is missing; what can be said about its possible value? Since η_n is the sum of those a_i for which there is a 1 in the last row of S_n , we have

$$\eta_n \leq \sum_{i=1}^n a_i. \quad (75)$$

From (74) and (75) we obtain the bounds

$$0 \leq \eta_n \leq \frac{2}{n-1} \sum_{i=1}^{n-1} \eta_i. \quad (76)$$

A similar result can be given if two measurements are missing. In some cases the lower bound in (76) can be improved by considering the differences $\eta_i - \eta_j$ for suitably chosen i and j .

An analysis of this type can also be given for other multiplexing schemes, including Fourier spectroscopy. A single missing measurement never entails complete loss of information about the spectrum. The uncertainty however increases rapidly as the number of missing data points increases.

6.3 Correction Procedures

A similar situation arises when one or more measurements are lost because of a burst of noise. A single noise spike of amplitude δ produces random noise over the whole spectrum with r.m.s. value $2\delta/n$, as we saw in §5.

A large noise spike can often be recognized by examining the other measurements obtained in the same run. When the spectrum is continuous, or contains a large number of intense lines, the individual data values do not greatly deviate from each other. Only if one or two spectral lines dominate do wide deviations occur. For a polychromatic spectrum, then, large noise spikes tend to appear as well defined, extraordinary points.

In the laboratory we have tended to remove these spikes and replace them by the mean value of adjacent data points on either side. The theoretical justification for this procedure is somewhat questionable, but we find in practice that the spectra obtained by means of such corrections are rather good approximations to the expected forms. The procedure is illustrated in Figs. 9-12. Figure 9 shows a spectrum obtained in the usual way with a 255×255 S-matrix. Figure 10 shows the effect of a noise spike in the last measurement, while Fig. 11 shows the result of replacing this and four adjacent data points by a straight line joining the 250-th and 1-st data points. Similarly in Fig. 12 the last 15 measurements have been replaced by 15 points linearly interpolated between the 240-th and 1-st measurements.

A comparison of Figs. 9 and 11 suggests that linear interpolation between values spanning the gap of missing data entries produces relatively little distortion. A somewhat better procedure might be to use a quadratic or higher degree curve for the interpolation.

Figure 13 uses the same data as in Fig. 9, except that a correction has been inserted to remove the negative echo caused by faulty slit deposition (or etching), as in Tai et al. (1975).

Acknowledgements

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Figure Captions

Fig. 1. Impulse response $H(\nu - \nu_0)$ (a) when diffraction can be ignored, and (b) when diffraction dominates. (c) and (d) are the corresponding slit functions. (See Examples (1) and (2) in §2.4.)

Fig. 2. A piecewise linear input spectrum $F(\nu)$. a_i is the intensity at the midpoint of the interval θ, θ_{i+1} .

Fig. 3. Showing how a sharp input (a) is coarsened (b) if the correction matrix T^{-1} is not used.

Fig. 4. Faulty mask alignment.

Fig. 5. Computer simulation of the errors produced by a systematic difference between slit width and step size. The final position of the mask after 255 steps is displaced by one slit width from the correct position. An input spectrum with a single sharp spectral line has been broadened to a width of two spectral elements, and noise has appeared across the whole spectrum.

Fig. 6. Mask designed to be moved continuously. This is a mask of length 14 corresponding to the W matrix of Eq. (50). Elements 1-6 are repeated as elements 8-13 respectively. The extra half slit at each end, corresponding to element 7, is needed to avoid errors at the ends of the spectrum.

Fig. 7. Path swept out by an open slit crossing the j -th segment θ, θ_{j+1} during the i -th detector reading. Time runs downward in this diagram, and the mask moves to the right.

Fig. 8. The ideal case when the mask has the correct dimensions, is properly aligned, and moves at the right speed. Again time runs downward and the mask moves to the right.

Fig. 9. Uncorrected 255-element spectrum of a laboratory source showing the 1.7μ mercury vapor lines. There has been no data processing other than applying the inverse Hadamard transformation.

Fig. 10. Same as Fig. 9 except that the last data point, of height 1.16×10^3 , has been replaced by a noise spike of height 4×10^3 . The main spectral features remain, but noise is added to the entire spectrum.

Fig. 11. Same as Fig. 9, except that we have simulated losing the last five data points and replacing them by a straight line joining the 250-th and 1-st data points.

Fig. 12. Same as Fig. 11, but now the last 15 measurements have been replaced by 15 points linearly interpolated between the 240-th and 1-st data points.

Fig. 13. The data used here is the same as in Fig. 9, except that a correction has been inserted to remove the negative echo caused by faulty slit deposition (or etching), as in Tai et al. (1975).

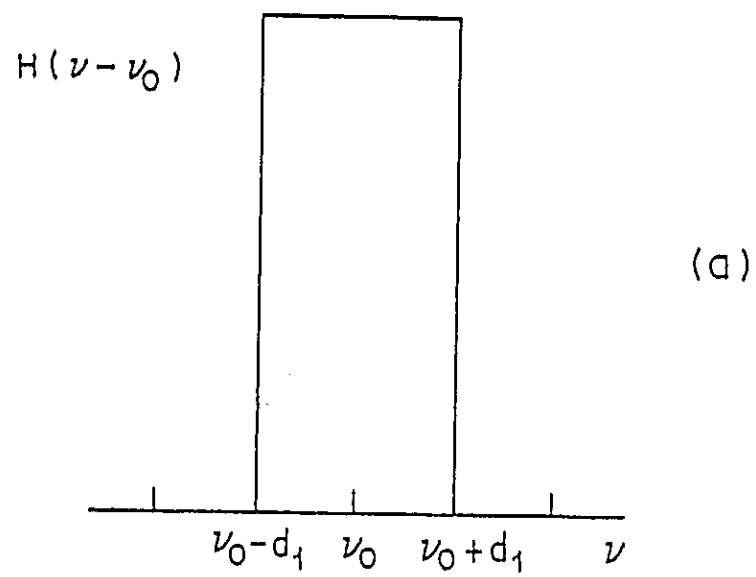
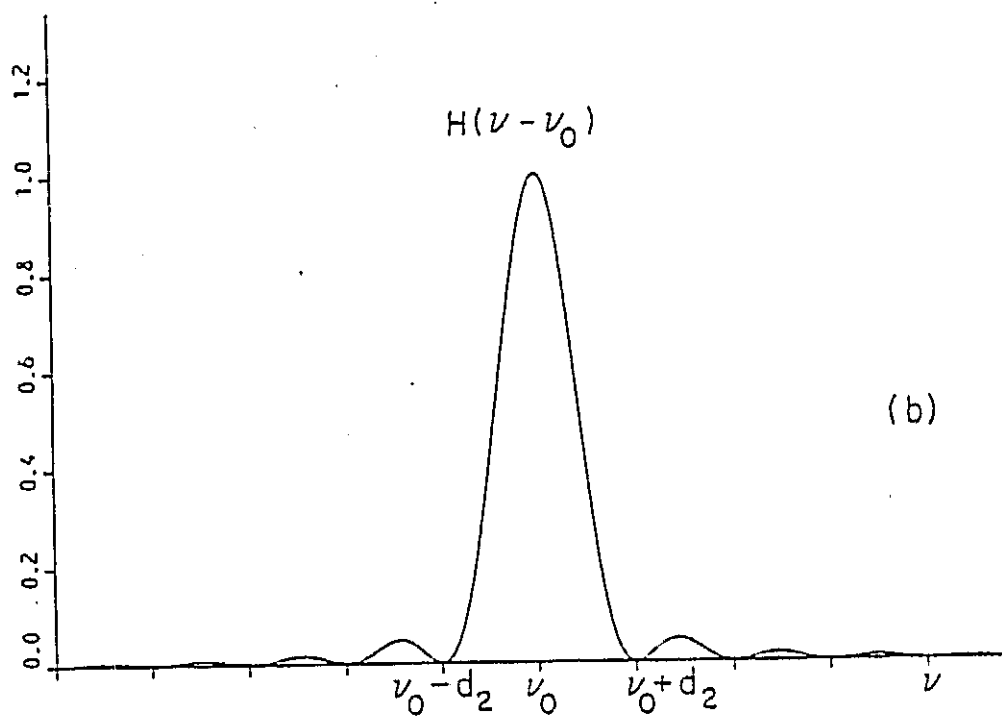


FIG 1(a)

FIG. 1(b)



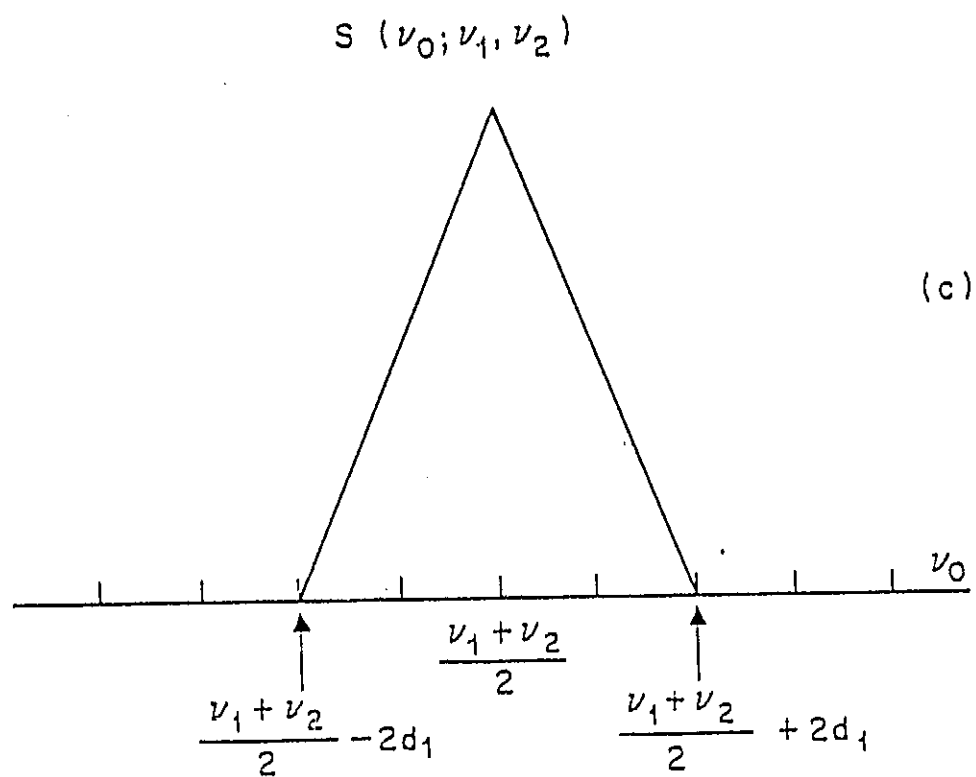
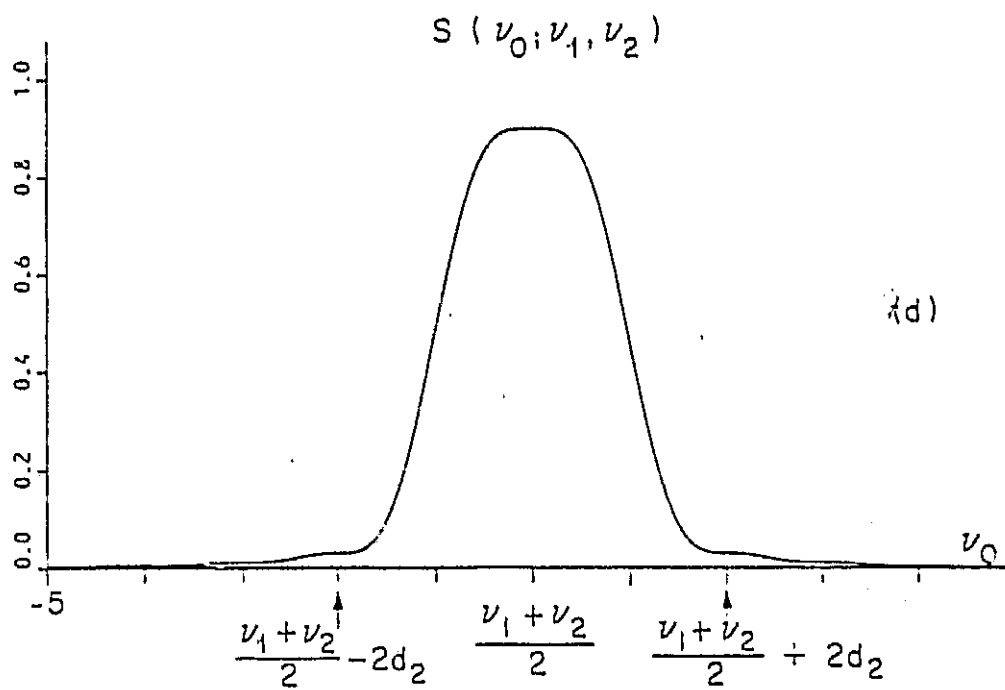


FIG 1(c)

FIG. 1(d)



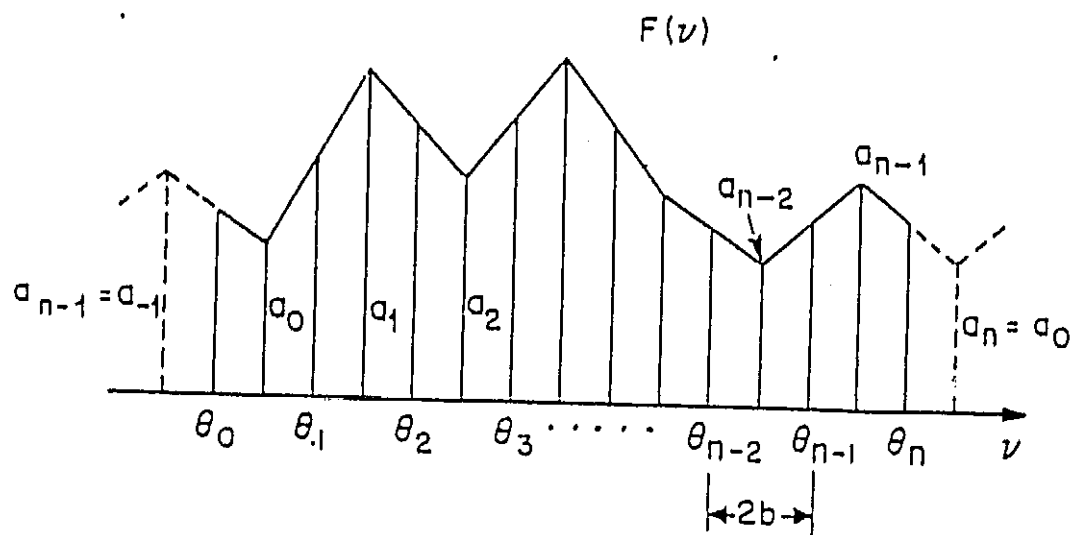


FIG. 2

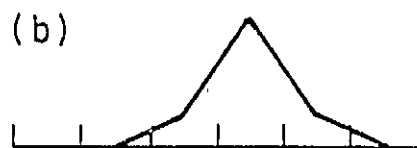
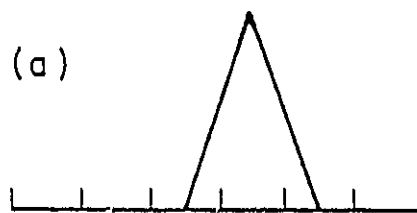
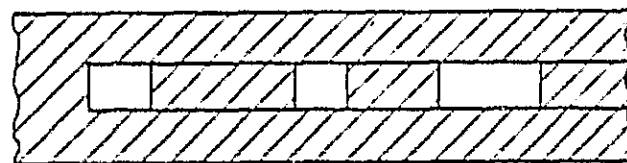
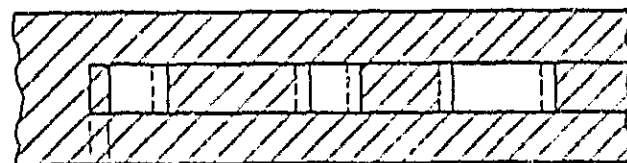


FIG. 3



CORRECT
POSITION

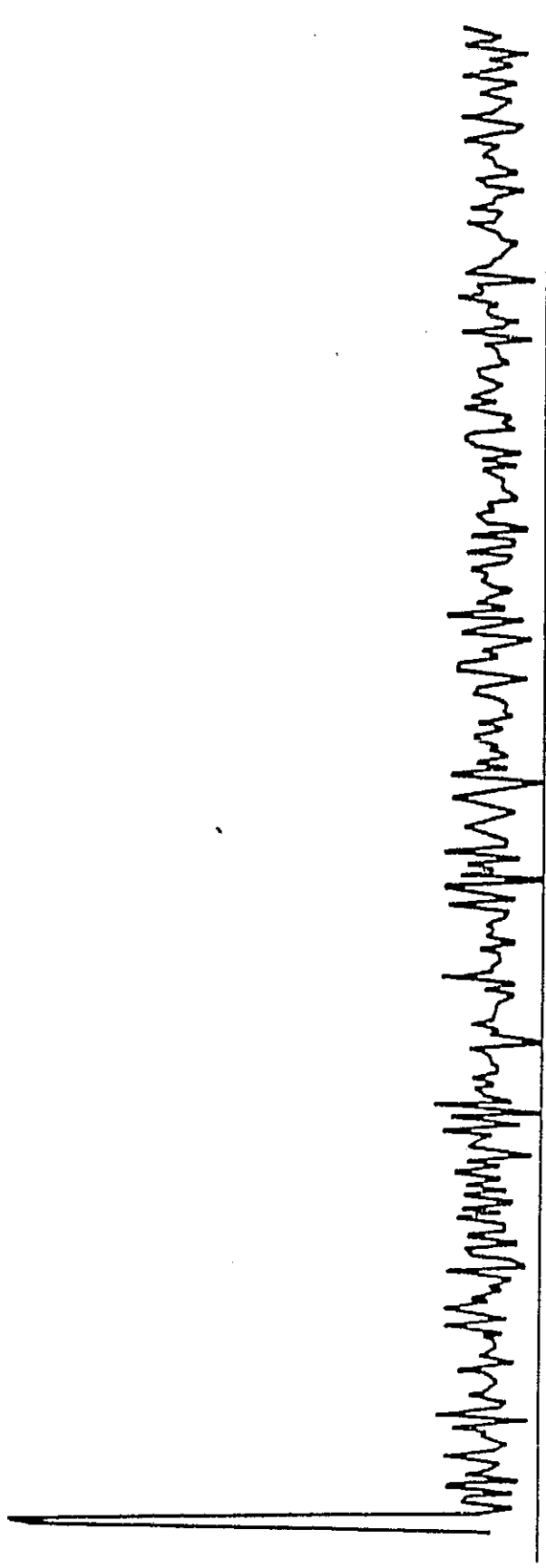


MISALIGNED



FIG. 4

FIG 5



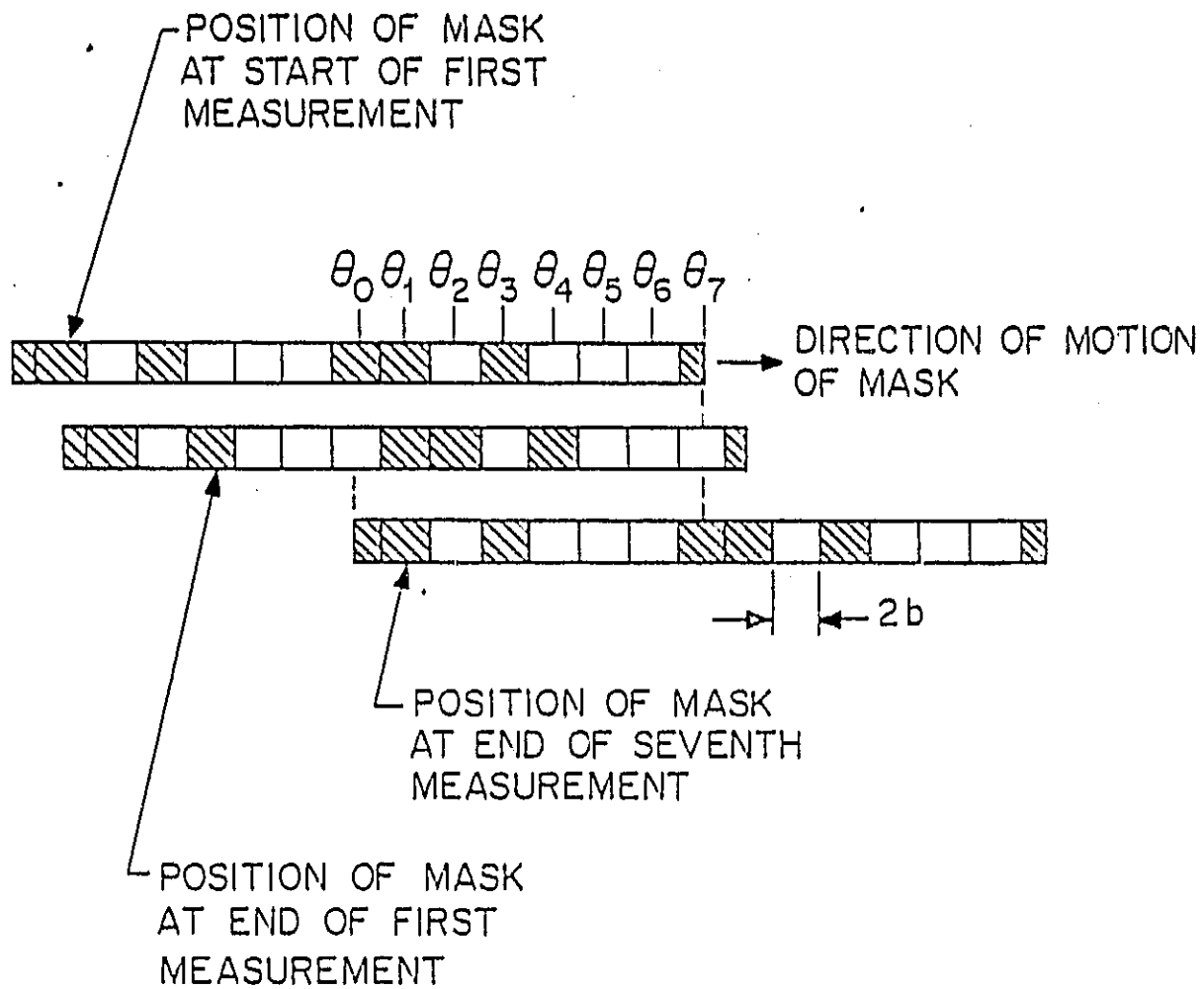


FIG. 6

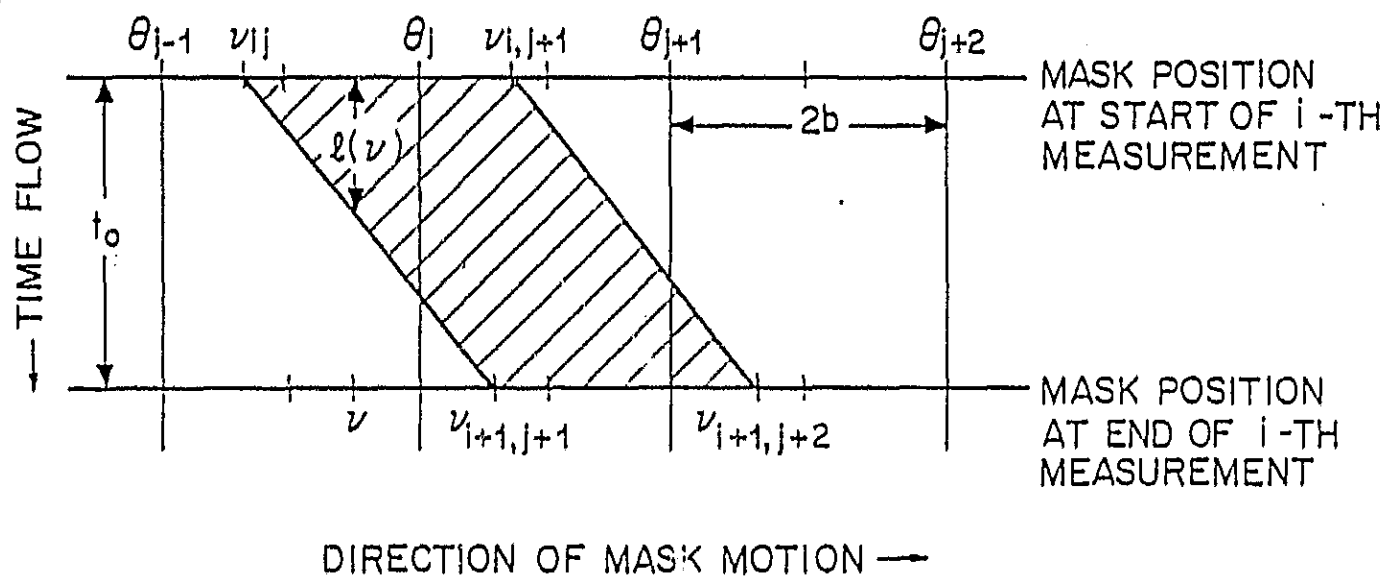


FIG. 7

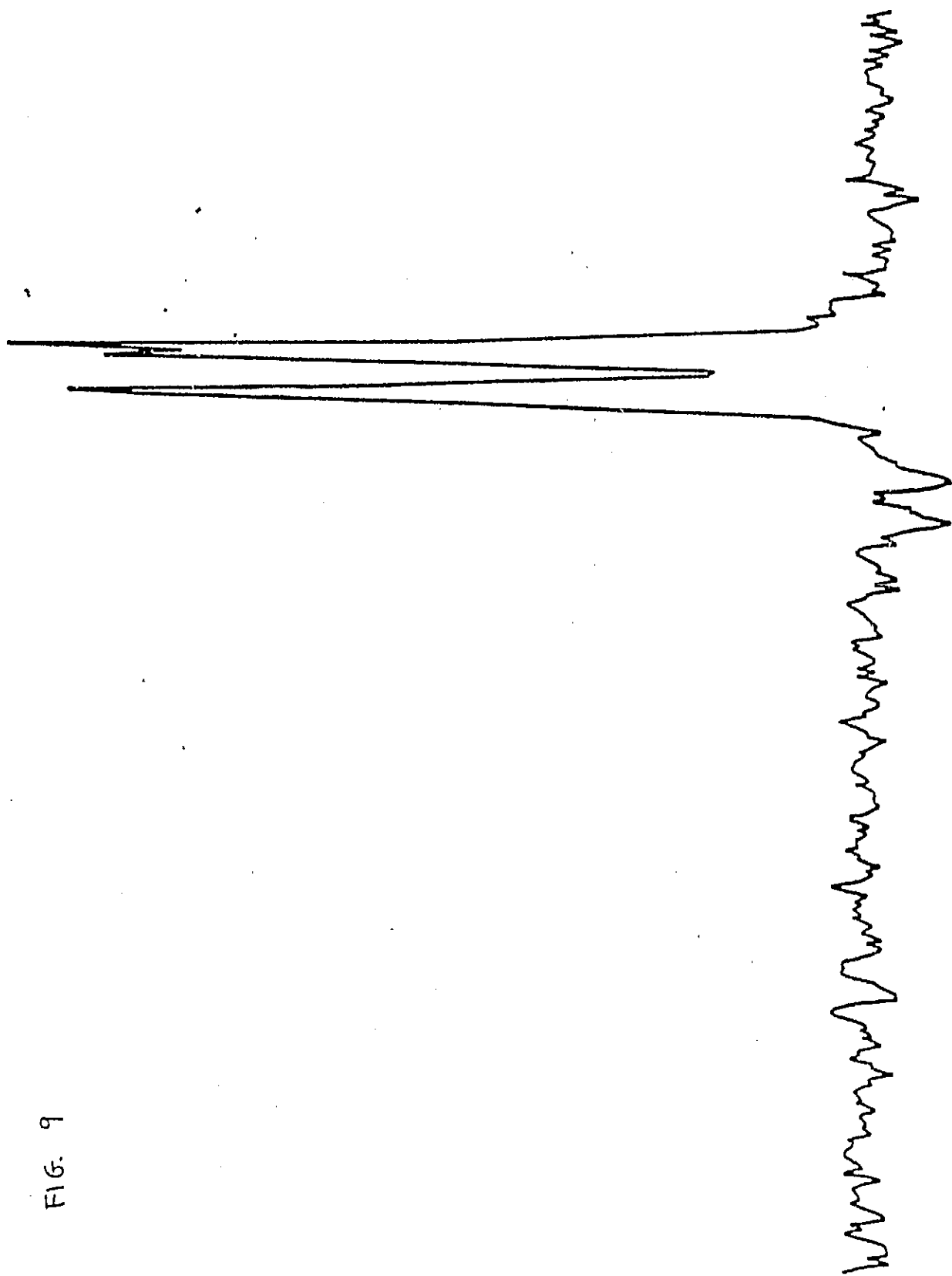


FIG. 9

FIG. 10

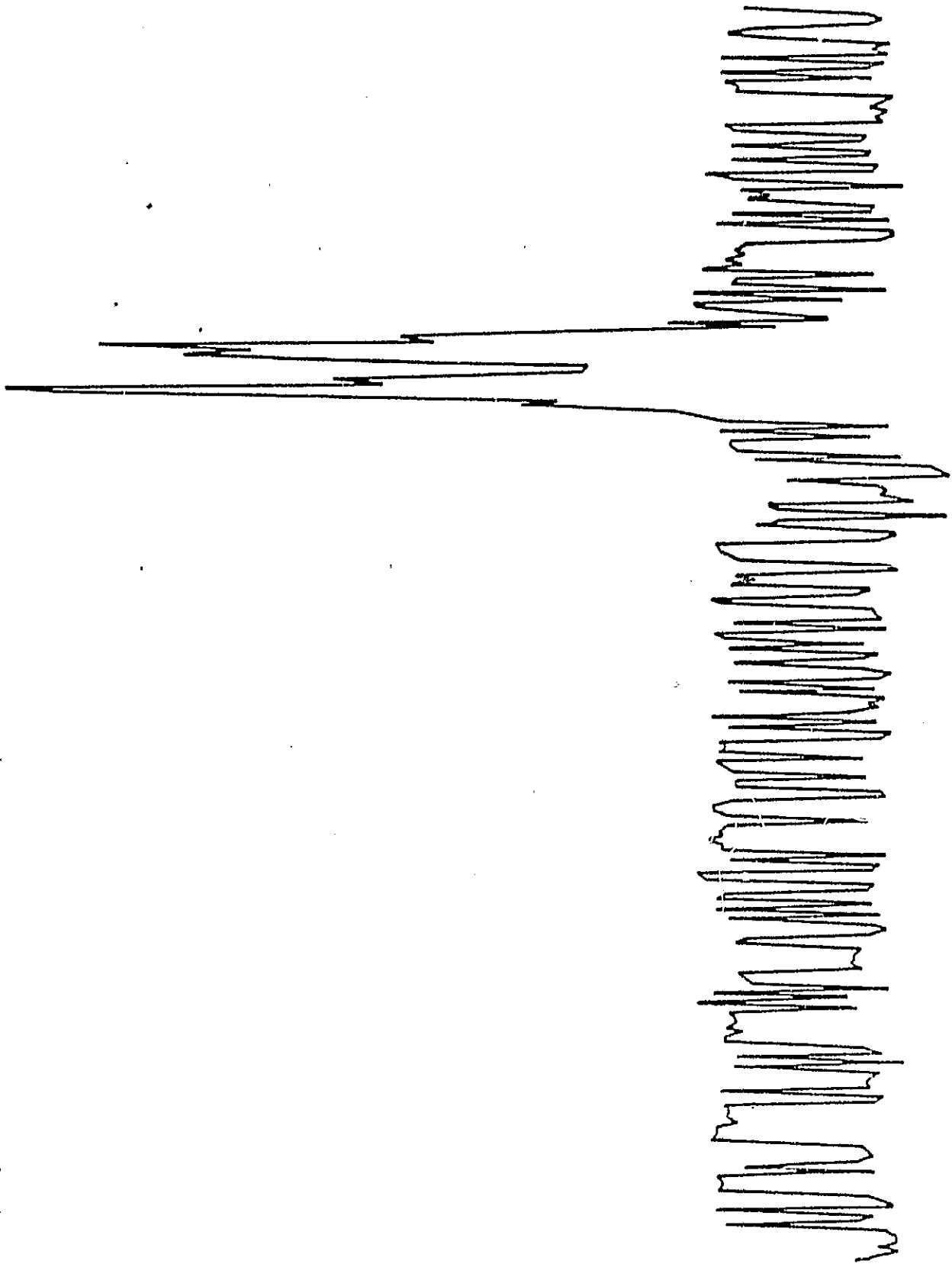


FIG. 11

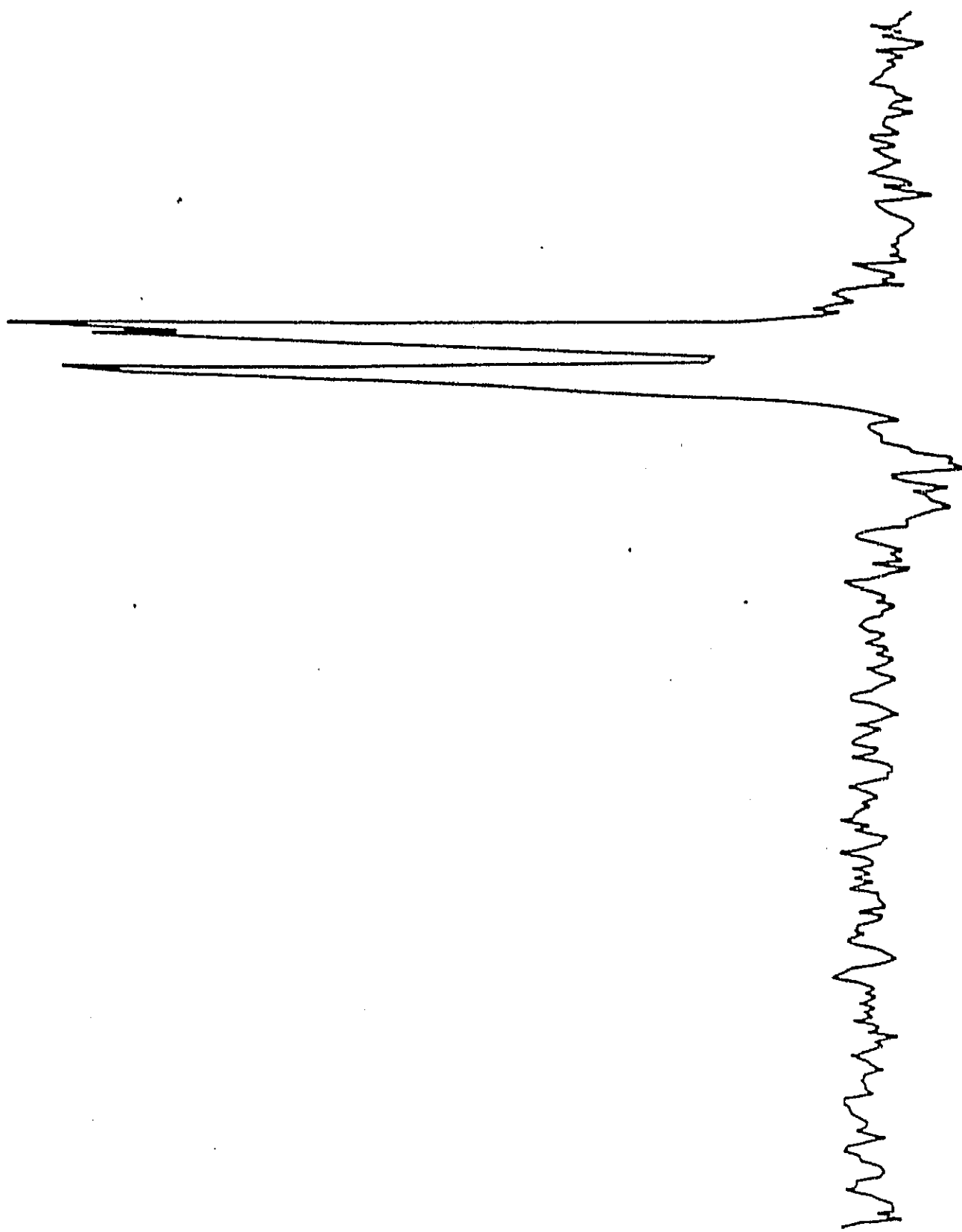


FIG. 12

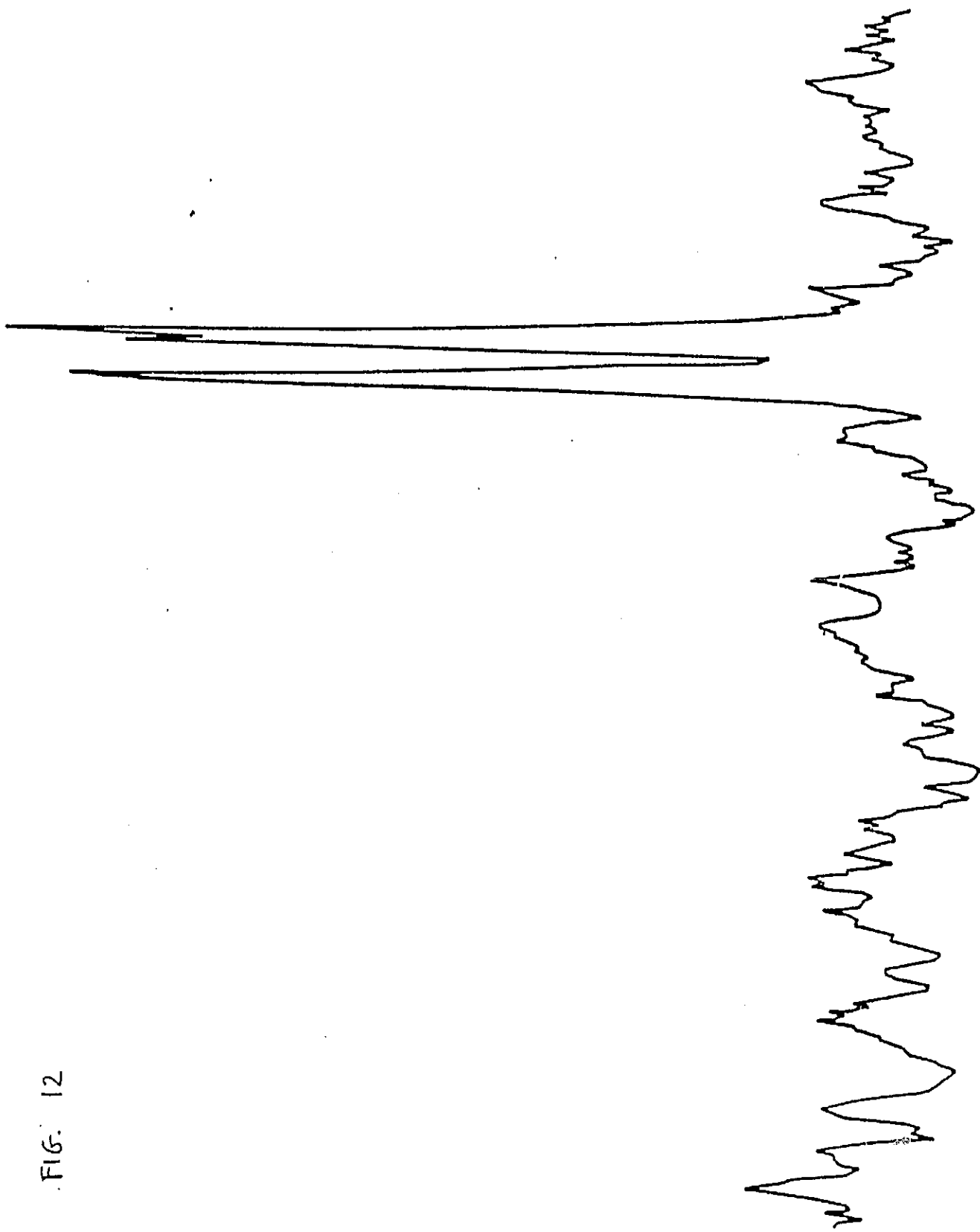
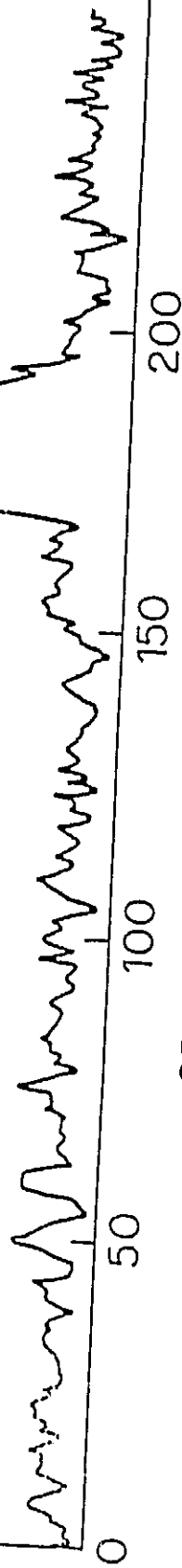


FIG. 13

INTENSITY

Hg VAPOR EMISSION LINE
SINGLE ENTRANCE, 255 EXIT SLITS

1.7 μ DOUBLET
EMISSION



SPECTRAL ELEMENT NUMBER

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